# Lipschitz stability in the lateral Cauchy problem for elasticity system 

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$$
\begin{aligned}
& \text { Abstract } \\
& \text { We consider the isotropic elasticity system: } \\
& \qquad \begin{aligned}
& \rho \partial_{t}^{2} \mathbf{u}-\mu\left(\Delta \mathbf{u}+\nabla\left(\nabla^{T} \mathbf{u}\right)\right)-\nabla\left(\lambda \nabla^{T} \mathbf{u}\right) \\
&-\sum_{j=1}^{3} \nabla \mu \cdot\left(\nabla u_{j}+\partial_{j} \mathbf{u}\right) \mathbf{e}_{j}=0 \quad \text { in } \quad \Omega \times(0, T)
\end{aligned}
\end{aligned}
$$

for the displacement vector $\mathbf{u}=\left(u_{1}, u_{2}, u_{3}\right)$ depending on $x \in \Omega$ and $t \in(0, T)$ where $\Omega$ is a bounded domain in $\mathbb{R}^{3}$ with the $C^{2}$-boundary, and we assume the density $\rho \in C^{2}(\bar{\Omega} \times[0, T])$ and the Lamé parameters $\mu, \lambda \in C^{3}(\bar{\Omega} \times[0, T])$. We will give Lipschitz stability estimates for solutions $\mathbf{u}$ to the above elasticity system with the lateral boundary data

$$
\mathbf{u}=\mathbf{g} \quad \text { on } \partial \Omega \times(0, T), \quad \partial_{\nu} \mathbf{u}=\mathbf{h} \quad \text { on } \Gamma \times(0, T)
$$

where $\Gamma$ is some part of $\partial \Omega$. Our proof is based on (1) a Carleman estimate with boundary data, (2) cut-off technique, and (3) principal diagonalization of the Lamé system.

## 1. Introduction and basic results

In this paper we are interested in Lipschitz stability estimates in the lateral Cauchy problem for the classical elasticity system

$$
\begin{align*}
\rho \partial_{t}^{2} \mathbf{u} & -\mu\left(\Delta \mathbf{u}+\nabla\left(\nabla^{T} \mathbf{u}\right)\right)-\nabla\left(\lambda \nabla^{T} \mathbf{u}\right) \\
& -\sum_{j=1}^{3} \nabla \mu \cdot\left(\nabla u_{j}+\partial_{j} \mathbf{u}\right) \mathbf{e}_{j}=0 \quad \text { in } \quad \Omega \times(0, T) \tag{1.1}
\end{align*}
$$

for the displacement vector $\mathbf{u}=\left(u_{1}, u_{2}, u_{3}\right)$ depending on $x=\left(x_{1}, x_{2}, x_{3}\right) \in \Omega$ and $t \in(0, T)$. Here $\Omega$ is a bounded domain in $\mathbb{R}^{3}$ with the $C^{2}$-boundary, $Q \equiv \Omega \times(0, T), \nu$ is the outward unit normal vector to $\partial \Omega$, and $\mathbf{e}_{1}=\left(\begin{array}{l}1 \\ 0 \\ 0\end{array}\right)$,
$\mathbf{e}_{2}=\left(\begin{array}{l}0 \\ 1 \\ 0\end{array}\right), \mathbf{e}_{3}=\left(\begin{array}{l}0 \\ 0 \\ 1\end{array}\right)$. We will assume that the density $\rho \in C^{2}(\bar{Q})$ and the Lamé parameters $\mu, \lambda \in C^{3}(\bar{Q}), \mu>0$ and $3 \lambda+2 \mu>0, \rho>0$ on $\bar{Q}$. We will give stability estimates for solutions $\mathbf{u}$ to (1.1) with the lateral boundary data

$$
\begin{equation*}
\mathbf{u}=\mathbf{g} \quad \text { on } \partial \Omega \times(0, T), \quad \partial_{\nu} \mathbf{u}=\mathbf{h} \quad \text { on } \Gamma \times(0, T), \tag{1.2}
\end{equation*}
$$

where $\Gamma$ is some part of $\partial \Omega$.
We will use the following notation: $\partial_{j}=\partial / \partial x_{j}, \partial_{t}=\partial / \partial t, \nabla=\left(\partial_{1}, \ldots\right.$, $\left.\partial_{n}\right), \nabla_{x, t}=\left(\partial_{1}, \ldots, \partial_{n}, \partial_{t}\right), \square_{a}=a \partial_{t}^{2}-\Delta, \Sigma=\partial \Omega \times(0, T)$. We set

$$
\|u\|_{2}(\Omega)=\left(\int_{\Omega} u^{2} d x\right)^{1 / 2}
$$

$H^{s}(\Omega), s \in \mathbb{N}$ denotes the $L^{2}$-based Sobolev space with the norm

$$
\|u\|_{(s)}^{2}(\Omega)=\sum_{|\alpha| \leq s}\left\|D^{\alpha} u\right\|_{2}^{2}(\Omega) .
$$

When we consider the space $H^{s}\left(\mathbb{R}^{n}\right)$, this norm can be equivalently defined by means of

$$
\|u\|_{(s)}^{2}\left(\mathbb{R}^{n}\right)=(2 \pi)^{-n} \int\left(1+|\xi|^{2}\right)^{s}|\widehat{u}(\xi)|^{2} d \xi
$$

where $\widehat{u}$ is the Fourier transform of $u$. This definition extends to $s \in \mathbb{R}$. We will drop the symbol $\Omega$ in norms and integrals when $\Omega=\mathbb{R}^{n}$. The norms $\|\cdot\|_{2}(\Gamma \times(0, T)),\|\cdot\|_{(1)}(\Gamma \times(0, T))$, etc. are defined similarly.

Our main results are the following theorems.
Theorem 1.1. Assume that the functions $\varphi, \varphi_{1}$ are strongly pseudoconvex correspondingly in $\bar{Q} \backslash\left(x^{0}, T / 2\right)$ and in $\bar{Q} \backslash\left(x^{1}, T / 2\right)$, with respect to the differential operators $\rho \partial_{t}^{2}-\mu \Delta, \rho \partial_{t}^{2}-(2 \mu+\lambda) \Delta, x^{0}, x^{1}$ are different points of $\Omega, \varphi_{1}\left(x^{1}, T / 2\right)<\varphi_{1}\left(x^{0}, T / 2\right)$,

$$
\begin{equation*}
\varphi<0, \quad \varphi_{1}<0 \quad \text { on } \quad \bar{\Omega} \times\{0, T\} \tag{1.3}
\end{equation*}
$$

and

$$
\begin{equation*}
0<\varphi, \quad 0<\varphi_{1} \quad \text { on } \quad \bar{\Omega} \times\{T / 2\} . \tag{1.4}
\end{equation*}
$$

Assume that $\Gamma \subset \partial \Omega$ satisfies

$$
\begin{align*}
\{x \in \partial \Omega ; & (\nabla \varphi(x, t) \cdot \nu(x))<0,  \tag{1.5}\\
& \left.\left(\nabla \varphi_{1}(x, t) \cdot \nu(x)\right)<0 \quad \text { for all } t \in[0, T]\right\} \supset \overline{\partial \Omega \backslash \Gamma} .
\end{align*}
$$

Then there is a constant $C$ depending only on $\Omega, \varphi, \varphi_{1}, \rho, \mu, \lambda, \Gamma$ such that

$$
\begin{align*}
& \left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} \nabla^{T} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\left\|\nabla^{T} \mathbf{u}(\cdot, t)\right\|_{(1)}(\Omega)  \tag{1.6}\\
& \quad \leq C\left\{\|\mathbf{h}\|_{2}(\Gamma \times(0, T))+\|\mathbf{g}\|_{(1)}(\Sigma)\right. \\
& \left.\quad+\left\|\nabla^{T} \mathbf{u}\right\|_{(1)}(\Sigma)+\left\|\partial_{\nu} \nabla^{T} \mathbf{u}\right\|_{2}(\Gamma \times(0, T))\right\}
\end{align*}
$$

when $0<t<T$.
If $x_{0} \notin \bar{\Omega}$, then we can take $\varphi_{1} \equiv \varphi$ and in Section 4 , we consider only such a case.

In the case of $x_{0} \in \bar{\Omega}$, a usual choice of $\varphi$ (e.g., see (4.3)) does not satisfy the pseudo-convexity. Therefore we need two weight functions $\varphi$ and $\varphi_{1}$ in Theorem 1.1. We will recall the definition of a strongly pseudo-convex function in Section 2.

We can give more explicit results in two important particular cases:
(1) $\Omega \subset B(0 ; R) \equiv\{x ;|x|<R\}$ and $\Gamma=\partial \Omega$,
(2) $\Omega \subset\left\{-h<x_{3}<0, x_{1}^{2}+x_{2}^{2}<r^{2}\right\}$ and $\Gamma$ is open in $\partial \Omega$.

Theorem 1.2. Assume that the functions $a=\rho / \mu$ and $a=\rho /(\lambda+2 \mu)$ satisfy the conditions

$$
\begin{align*}
& \theta^{2} a\left(a+\frac{1}{2}\left(t-\frac{T}{2}\right) \partial_{t} a+a^{-1 / 2}\left|\left(t-\frac{T}{2}\right) \nabla a\right|\right)<a+\frac{1}{2} x \cdot \nabla a-\frac{\beta}{2} \partial_{n} a  \tag{1.7}\\
& \theta^{2} a \leq 1 \text { on } \bar{Q} .
\end{align*}
$$

If in case (1) $R<\theta T / 2$ and in case (2) $\Gamma$ contains $\left\{0 \leq x \cdot \nu-\beta \nu_{n}\right\} \cap \partial \Omega$, $r^{2}+h^{2}+2 h \beta<\theta^{2} T^{2} / 4$, then bound (1.6) holds.

We expect that the terms with $\nabla^{T} \mathbf{u}$ are not needed in (1.6) and can be eliminated by more sophisticated methods.

Next in place of $\partial_{\nu} \mathbf{u}=\mathbf{h}$ in (1.2), we take the surface traction for a stability estimate. In this paper, for brevity, we further assume that $\mathbf{g}=0$. That is, in place of (1.2), we consider

$$
\begin{equation*}
\mathbf{u}=0 \quad \text { on } \partial \Omega \times(0, T), \quad \sigma(\mathbf{u}) \nu=\mathbf{h} \quad \text { on } \Gamma \times(0, T) . \tag{1.2}
\end{equation*}
$$

Here the stress tensor $\sigma(\mathbf{u})$, a $3 \times 3$ matrix, is given by

$$
\begin{aligned}
\sigma(\mathbf{u})= & \lambda\left(\nabla^{T} \mathbf{u}\right)\left(\mathbf{e}_{1}, \mathbf{e}_{2}, \mathbf{e}_{3}\right) \\
& +\mu\left(\nabla u_{1}+\partial_{1} \mathbf{u}, \nabla u_{2}+\partial_{2} \mathbf{u}, \nabla u_{3}+\partial_{3} \mathbf{u}\right) .
\end{aligned}
$$

Theorem 1.3. Assume that the functions $\varphi, \varphi_{1}$ are strongly pseudoconvex with respect to the differential operators $\rho \partial_{t}^{2}-\mu \Delta, \rho \partial_{t}^{2}-(2 \mu+\lambda) \Delta$ correspondingly on $\bar{Q} \backslash\left(x^{0}, T / 2\right)$, and on $\bar{Q} \backslash\left(x^{1}, T / 2\right), x^{0}, x^{1}$ are two different points of $\Omega, \varphi_{1}\left(x^{1}, T / 2\right)<\varphi_{1}\left(x^{0}, T / 2\right)$, and satisfy (1.3), (1.4), and (1.5). Then there is a constant $C$ depending only on $\Omega, \varphi, \varphi_{1}, \rho, \mu, \lambda, \Gamma$ such that

$$
\begin{align*}
&\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} \nabla^{T} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega) \\
& \quad+\left\|\nabla^{T} \mathbf{u}(\cdot, t)\right\|_{(1)}(\Omega)  \tag{1.8}\\
& \leq C\left\{\|\mathbf{h}\|_{2}(\Gamma \times(0, T))+\|\nabla \mathbf{h}\|_{2}(\Gamma \times(0, T))\right\}
\end{align*}
$$

if $\mathbf{u} \in C^{2}(\bar{Q})^{3}$ satisfies (1.1) and (1.2)', and $\nabla^{T} \mathbf{u} \in C^{2}(\bar{Q})^{3}$.

Theorem 1.4. Assume that the functions $a=\rho / \mu$ and $a=\rho /(\lambda+2 \mu)$ satisfy conditions (1.7). If in case (1) $R<\theta T / 2$ and $\Gamma=\partial \Omega$, and in case (2) $\Gamma$ contains $\left\{0 \leq x \cdot \nu-\beta \nu_{n}\right\} \cap \partial \Omega, r^{2}+h^{2}+2 h \beta<\theta^{2} T^{2} / 4$, then bound (1.8) holds.

Sharp uniqueness in the continuation for the static elasticity system with any first order perturbations (and smooth coefficients) was obtained by Dehman and Robbiano [6]. Using much more elementary techniques, Ang, Ikehata, Trong and Yamamoto ([2]) reduced a static system with zero order perturbations to a principally diagonal one and obtained sharp uniqueness of the continuation results under reduced smoothness assumptions. The dynamical system with constant coefficients has been considered by Bukhgeim and Kardakov [4] by using spherical means and by Alabau and Komornik [1] who transferred to this system the original approach of multipliers used by Lop Fat Ho [22] for the wave equation to obtain sharp stability estimates in energy norms (cf. Grasselli and Yamamoto [8]). The time dependent classical elasticity with variable coefficients and arbitrary first order perturbations was considered by Isakov [13] by pseudo-convexity methods. See also Isakov [14]. In the forthcoming paper (Eller, Isakov, Nakamura and Tataru [7]), sharp uniqueness results for the lateral Cauchy problem for system (1.1) and some conditional Hölder type stability estimates are proved on the basis of Tataru [25]. The results of these papers are obtained by combining known results for scalar hyperbolic equations of second order and principal diagonalization of the elasticity system and they imply approximate boundary controllability of this system. For the stabilization for the elasticity system, we can refer to Horn [11]. Yamamoto [28] applied principal diagonalization to some inverse source problems for the Maxwell system. Recently, Lebeau and Zuazua [21] obtained quite complete interior controllability results for a thermoelasticity system where Lipschitz stability is not possible. About general approaches to such estimates for scalar hyperbolic equations of second order, we refer to papers of Bardos, Lebeau and Rauch [3], Burq [5] (methods of geometrical optics) and Tataru [26] (methods of Carleman estimates).

## 2. Lipschitz stability for principally diagonal hyperbolic systems

In this section, we do not assume that the spatial dimension is 3 . In other words, $\Omega$ is a bounded domain in an $n$-dimensional space.

We are considering the following system of partial differential equations

$$
\begin{equation*}
P(j) u_{j}+b_{j}\left(x, t ; \nabla_{x, t} \mathbf{u}, \mathbf{u}\right)=0, \quad j=1, \ldots, m \quad \text { in } Q=\Omega \times(0, T), \tag{2.1}
\end{equation*}
$$

where the principal part $P(j)$ is a real second order $t$-hyperbolic operator with $C^{1}$-smooth coefficients, $b_{j}$ are linear functions of $\nabla \mathbf{u}$ with $L^{\infty}(Q)$-coefficients and $\mathbf{u}$ with $L^{n+1}(Q)$-coefficient, and $\mathbf{u}=\left(u_{1}, \ldots, u_{m}\right)$. We will impose the following Dirichlet lateral boundary condition

$$
\begin{equation*}
\mathbf{u}=\mathbf{g} \quad \text { on } \quad \Sigma . \tag{2.2}
\end{equation*}
$$

To formulate results we need a (weight) function $\varphi \in C^{2}(\bar{Q})$. This function is called strongly pseudo-convex on $\bar{Q}$ with respect to the operator $P$ with the principal symbol $p(x, t ; \zeta), \zeta \in \mathbb{C}^{n+1}$ if $\nabla \varphi$ is not zero on $\bar{Q}$ and if for any $(x, t) \in \bar{Q}$, the equality

$$
p(x, t ; \zeta)=0, \quad \zeta=\xi+i \tau \nabla \varphi, \quad \tau \neq 0
$$

or

$$
p(x, t ; \xi)=0, \quad \nabla_{\xi} p(x, t ; \xi) \cdot \nabla_{x, t} \varphi(x, t)=0, \quad \xi \neq 0
$$

implies that

$$
\sum_{j, k=1}^{n+1} \partial_{j} \partial_{k} \varphi \partial p / \partial \zeta_{j} \overline{\partial p / \partial \zeta_{k}}+\tau^{-1} \Im\left(\partial_{k} p \overline{\partial p / \partial \zeta_{k}}\right)>0
$$

(cf. Hörmander [9], [10]). Here $\bar{\alpha}$ denotes the complex conjugate and $\Im$ is the imaginary part.

Let $\Gamma_{-}=\{x \in \partial \Omega:(\nabla \varphi(x, t) \cdot \nu(x))<0$ for all $t \in(0, T)\}, \Gamma_{+}=\partial \Omega \backslash \Gamma_{-}$ and $\Gamma$ be a neighbourhood of $\Gamma_{+}$in $\partial \Omega$.

Theorem 2.1. Assume that $\varphi, \varphi_{1}$ are strongly pseudo-convex correspondingly on $\bar{Q} \backslash\left(x^{0}, T / 2\right)$ and on $\bar{Q} \backslash\left(x^{1}, T / 2\right)$ with respect to $P(j), j=$ $1, \ldots, m$, and satisfy conditions (1.3), (1.4), and (1.5), $x^{0}$, $x^{1}$ are two different points of $\Omega, \varphi_{1}\left(x^{1}, T / 2\right)<\varphi_{1}\left(x^{0}, T / 2\right)$. Let $\mathbf{u}$ satisfy (2.1) and (2.2) with $\mathbf{g}=0$. Then there is a constant $C$ depending only on $\Omega, \varphi, \varphi_{1}, P(j), b_{j}, \Gamma$ such that

$$
\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega) \leq C\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Gamma \times(0, T))
$$

for all $t \in(0, T)$.
Proof. By conditions (1.3) and (1.4), there is $\varepsilon_{1}>0$ such that

$$
\begin{equation*}
\varphi<-\varepsilon_{1}, \quad \varphi_{1}<-\varepsilon_{1} \quad \text { on } \quad \Omega \times\left(\left(0, \varepsilon_{1}\right) \cup\left(T-\varepsilon_{1}, T\right)\right) \tag{2.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\varepsilon_{1}<\varphi, \quad \varepsilon_{1}<\varphi_{1} \quad \text { on } \quad \Omega \times\left(T / 2-\varepsilon_{1}, T / 2+\varepsilon_{1}\right) \tag{2.4}
\end{equation*}
$$

Henceforth we set $B\left(x_{0}, T / 2 ; \delta\right)=\left\{(x, t) \in \mathbb{R}^{n+1} ;\left|x-x_{0}\right|^{2}+|t-T / 2|^{2}<\right.$ $\left.\delta^{2}\right\}$. To apply Tataru's Carleman estimates ([27, Theorem 1]), we make use of cut-off functions $0 \leq \chi, \chi_{1}, \chi_{2} \leq 1$. Let $\chi \in C^{\infty}(\mathbb{R}), \chi=0$ on $\left(0, \varepsilon_{1} / 2\right) \cup$ $\left(T-\varepsilon_{1} / 2, T\right), \chi=1$ on $\left(\varepsilon_{1}, T-\varepsilon_{1}\right)$. Let $\chi_{1} \in C^{\infty}\left(\mathbb{R}^{n}\right), \chi_{1}=1$ near $\partial \Omega \backslash \Gamma$ and $\chi_{1}=0$ near $\Gamma_{+}$. Let $B_{1}$ be the ball $B\left(x^{0}, T / 2 ; \delta\right) \subset Q$ with $\bar{B}_{1}$ not containing the point $\left(x_{1}, T / 2\right)$. Let $\chi_{2} \in C^{\infty}\left(\mathbb{R}^{n+1}\right), \chi_{2}=1$ on $Q \backslash B_{1}, \chi_{2}=0$ on $B_{0} \equiv B\left(x^{0}, T / 2 ; \delta / 2\right)$. In addition using the conditions of Theorem 2.1, we can choose $\delta$ so small that $\varphi_{1}\left(x^{1}, T / 2\right)<\varphi_{1}(x, t)$ for $(x, t) \in B_{1}$.

We have

$$
\begin{gather*}
\mathbf{u}=\mathbf{u}_{0}+\mathbf{u}_{1}+\mathbf{u}_{2}+\mathbf{u}_{3}, \quad \text { where } \\
\mathbf{u}_{0}=(1-\chi) \chi_{2} \mathbf{u}, \quad \mathbf{u}_{1}=\chi \chi_{1} \chi_{2} \mathbf{u}  \tag{2.5}\\
\mathbf{u}_{2}=\chi\left(1-\chi_{1}\right) \chi_{2} \mathbf{u}, \quad \mathbf{u}_{3}=\left(1-\chi_{2}\right) \mathbf{u} .
\end{gather*}
$$

We set $\mathbf{u}_{i}=\left(u_{i, 1}, \ldots, u_{i ; m}\right), i=0,1,2$.
Observe that supp $\mathbf{u}_{1} \subset K_{1}$ where $K_{1}$ is a compact subset of $\bar{\Omega} \times\left[\varepsilon_{1} / 2, T-\right.$ $\left.\varepsilon_{1} / 2\right] \backslash B_{0}$ with $K_{1} \cap \partial Q$ contained in $\Gamma_{2} \times[0, T] \subset \Gamma_{-} \times[0, T]$ which does not intersect $\bar{\Gamma}_{+}$. Since $\Gamma_{2} \subset \Gamma_{-}$and does not intersect $\bar{\Gamma}_{+}$, the Dirichlet boundary operator on $\Gamma_{2}$ satisfies the strong Lopatinskii condition in the direction $\nabla_{x, t} \varphi$ ([27, Proposition 5.1]). Henceforth let $|\alpha| \leq 1$. By Theorem 1 in [27], there exists a constant $C>0$ such that

$$
\begin{equation*}
\tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{1 ; j}\right\|_{2}^{2}(Q) \leq C\left\|e^{\tau \varphi} P(j) u_{1 ; j}\right\|_{2}^{2}(Q), \quad 1 \leq j \leq m \tag{2.6}
\end{equation*}
$$

for large $\tau$.
Similarly, $\operatorname{supp} \mathbf{u}_{2} \subset K_{2}$ where $K_{2}$ is a compact subset of $\bar{\Omega} \times\left[\varepsilon_{1} / 2, T-\right.$ $\left.\varepsilon_{1} / 2\right] \backslash B_{0}$ with $K_{2} \cap \partial Q$ contained in $\Gamma \times[0, T]$. The Cauchy boundary operators always satisfy the strong Lopatinskii condition (cf. [27]) in the direction $\nabla_{x, t} \varphi$. By Theorem 1 in [27], we have

$$
\begin{align*}
& \tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{2 ; j}\right\|_{2}^{2}(Q)  \tag{2.7}\\
& \quad \leq C\left\{\left\|e^{\tau \varphi} P(j) u_{2 ; j}\right\|_{2}^{2}(Q)+\tau\left\|e^{\tau \varphi} \partial_{\nu} u_{2 ; j}\right\|_{2}^{2}(\Gamma \times(0, T))\right\}, \quad 1 \leq j \leq m
\end{align*}
$$

for large $\tau>0$.
Here and henceforth $C>0$ denotes a generic constant which is independent of $\tau$. From (2.5) and Leibniz' formula

$$
P(j) u_{1 ; j}=\chi \chi_{1} \chi_{2} P(j) u_{j}+A_{1 ; j}\left(u_{j}\right) \equiv A_{1}(u ; j), \quad 1 \leq j \leq m,
$$

due to equations (2.1). Here $A_{1 ; j}, A(\cdot ; j)$ are linear partial differential operators of first order whose coefficients of the first-order terms are in $L^{\infty}(Q)$ and the coefficients of the zeroth-order terms are in $L^{n+1}(Q)$ (depending on $\chi, \chi_{1}$ and $\left.\chi_{2}\right)$. A similar formula holds for $u_{2 ; j}, 1 \leq j \leq m$.

Using these formulae, adding inequalities (2.6), (2.7) and summing over $j=1, \ldots, m$ and $|\alpha| \leq 1$, we arrive at

$$
\begin{align*}
& \sum_{j=1}^{m} \sum_{|\alpha| \leq 1} \tau^{3-2|\alpha|}\left\{\left\|e^{\tau \varphi} \partial^{\alpha} u_{1 ; j}\right\|_{2}^{2}(Q)+\left\|e^{\tau \varphi} \partial^{\alpha} u_{2 ; j}\right\|_{2}^{2}(Q)\right. \\
& \left.+\left\|e^{\tau \varphi} \partial^{\alpha} u_{0 ; j}\right\|_{2}^{2}(Q)+\left\|e^{\tau \varphi} \partial^{\alpha} u_{3 ; j}\right\|_{2}^{2}(Q)\right\} \\
& \leq
\end{align*} \quad C\left(\sum_{j=1}^{m} \sum_{|\alpha|=1}\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}(Q)+\sum_{j=1}^{m}\left\|q_{j} e^{\tau \varphi} u_{j}\right\|_{2}^{2}(Q) \text { ) } \quad \begin{array}{l}
\quad+\tau\left\|e^{\tau \varphi} \partial_{\nu} u_{j}\right\|_{2}^{2}(\Gamma \times(0, T))  \tag{2.8}\\
\left.\quad+\tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{0 ; j}\right\|_{2}^{2}(Q)+\left\|e^{\tau \varphi} \partial^{\alpha} u_{3 ; j}\right\|_{2}^{2}(Q)\right)
\end{array}\right.
$$

Here $q_{j} u_{j}$ is the zeroth-order term from $b_{1}, \ldots, b_{m}$. Let $n \geq 3$. Since $Q \in \mathbb{R}^{n+1}$ and $q_{j} \in L^{n+1}(Q)$ for $1=1, \ldots, m$, by Hölder's inequality and the Sobolev embedding we obtain

$$
\begin{aligned}
\left\|e^{\tau \varphi} q_{j} u_{j}\right\|_{2}^{2}(Q) & \leq\left\|q_{j}\right\|_{L^{n+1}(Q)}^{2}\left\|e^{\tau \varphi} u_{j}\right\|_{L^{2(n+1) /(n-1)}(Q)}^{2} \leq C\left\|e^{\tau \varphi} u_{j}\right\|_{H^{1}(Q)}^{2} \\
& \leq C\left(\tau^{2}\left\|e^{\tau \varphi} u_{j}\right\|_{2}^{2}(Q)+\sum_{|\alpha| \leq 1}\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}(Q)\right)
\end{aligned}
$$

Hence by taking $\tau$ large, we can eliminate the term $\left\|e^{\tau \varphi} q_{j} u_{j}\right\|_{2}^{2}(Q)$ in the right hand of (2.8), so that equality (2.5), the triangle inequality (in the left side of (2.8)), and (2.8) yield

$$
\begin{align*}
& \sum_{j=1}^{m} \sum_{|\alpha| \leq 1} \tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}(Q) \\
& \leq C\left(\sum _ { j = 1 } ^ { m } \sum _ { | \alpha | \leq 1 } \left\{\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}(Q)+\tau\left\|e^{\tau \varphi} \partial_{\nu} u_{j}\right\|_{2}^{2}(\Gamma \times(0, T))\right.\right.  \tag{2.9}\\
&\left.\left.+\tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}\left(Q(1) \cup B_{1}\right)\right\}\right)
\end{align*}
$$

where $Q(1)=\Omega \times\left(\left(0, \varepsilon_{1}\right) \cup\left(T-\varepsilon_{1}, T\right)\right)$. Choosing $\tau$ large, we can eliminate the first terms $\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}(Q)$ in the right side of (2.9).

Next we will eliminate the terms

$$
\tau^{3-2|\alpha|}\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}\left(\Omega \times\left(\left(0, \varepsilon_{1}\right) \cup\left(T-\varepsilon_{1}, T\right)\right)\right)
$$

Due to condition (2.4), the left side in (2.9) is greater than

$$
e^{2 \tau \varepsilon_{1}} \sum_{j=1}^{m} \sum_{|\alpha| \leq 1}\left\|\partial^{\alpha} u_{j}\right\|_{2}^{2}\left(\Omega \times\left(T / 2-\varepsilon_{1}, T / 2+\varepsilon_{1}\right)\right)
$$

Writing the last sum as

$$
\int_{T / 2-\varepsilon_{1}}^{T / 2+\varepsilon_{1}} E(t) d t, \quad E(t)=\sum_{j=1}^{m} \sum_{|\alpha| \leq 1}\left\|\partial^{\alpha} u_{j}(\cdot, t)\right\|_{2}^{2}(\Omega)
$$

and observing that by elementary properties, the integral with respect to $t$ is not less than $2 \varepsilon_{1} E(\theta)$ for some $\theta \in\left(T / 2-\varepsilon_{1}, T / 2+\varepsilon_{1}\right)$, we conclude that the left side in (2.9) is not less than $2 \varepsilon_{1} e^{2 \tau \varepsilon_{1}} E(\theta)$.

Since system (2.1) is $t$-hyperbolic, the known energy estimates (e.g., John [15]) imply that $E(t) \leq C E(\theta), E(t) \leq C E(0), 0 \leq t \leq T$. Using the above bounds from (2.3) and (2.9), we conclude that

$$
\begin{aligned}
e^{2 \tau \varepsilon_{1}} E(0) \leq & C\left(\tau\left\|e^{\tau \varphi} \partial_{\nu} \mathbf{u}\right\|_{2}^{2}(\Gamma \times(0, T))+\tau^{3} e^{-2 \tau \varepsilon_{1}} E(0)\right. \\
& \left.+\sum_{j=1}^{m} \sum_{|\alpha|=1} \tau\left\|e^{\tau \varphi} \partial^{\alpha} u_{j}\right\|_{2}^{2}\left(B_{1}\right)+\tau^{3}\left\|e^{\tau \varphi} \mathbf{u}\right\|_{2}^{2}\left(B_{1}\right)\right) .
\end{aligned}
$$

Again choosing $\tau$ large we can eliminate the term in the right side containing $E(0)$. After that we fix $\tau$ and, using that $E(t) \leq C E(0)$, we will have

$$
\begin{equation*}
E(t) \leq C\left(\left\|\partial_{\nu} \mathbf{u}\right\|_{2}^{2}(\Gamma \times(0, T))+\|\mathbf{u}\|_{(1)}^{2}\left(B_{1}\right)\right) . \tag{2.10}
\end{equation*}
$$

The last step of the proof is to eliminate the norm over $B_{1}$. To do it we will use the Hölder type conditional stability estimate in the Cauchy problem from [7].

Let $Q_{\varepsilon}=Q \cap\left\{(x, t) ; \varphi_{1}(x, t)-\varphi_{1}\left(x^{1}, T / 2\right)-\delta_{0}>\varepsilon\right\}$. Due to the choice of $B_{1}$ we have $B_{1} \subset Q_{\varepsilon}$ for some small positive $\varepsilon, \delta_{0}$ and $\left(x^{1}, T / 2\right)$ is not in $\bar{Q}_{0}$. We will fix such $\varepsilon, \delta_{0}$. By Theorem 3.3 in [7] we have

$$
\begin{aligned}
\|\mathbf{u}\|_{(1)}\left(B_{1}\right) & \leq\|\mathbf{u}\|_{(1)}\left(Q_{\varepsilon}\right) \\
& \leq C\left\{\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Gamma \times(0, T))+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}^{\theta}(\Gamma \times(0, T))\|\mathbf{u}\|_{(1)}^{1-\theta}\left(Q_{0}\right)\right\},
\end{aligned}
$$

where $\theta \in(0,1)$ depends on $\varepsilon$. Using the Hölder inequality, for any $\delta_{1}>0$, we can choose $C\left(\delta_{1}\right)>0$ such that $a^{\theta} b^{1-\theta} \leq \delta_{1} b+C\left(\delta_{1}\right) a$ and combining the above inequality with (2.10) without the norm over $Q(1)$ and with the obvious bound $\|\mathbf{u}\|_{(1)}^{2}(\Omega) \leq C E(0)$ we finally obtain

$$
\|\mathbf{u}\|_{(1)}(\Omega) \leq C\left(\delta_{1}\right)\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Gamma \times(0, T))+C \delta_{1}\|\mathbf{u}\|_{(1)}(Q)
$$

Choosing $\delta_{1}<1 / C$ we eliminate the last term and complete the proof of Theorem 2.1.

In this proof we have used general Carleman estimates of Tataru ([27]) and the device of Klibanov and Malinsky ([16]) and of Tataru ([26]) for deriving Lipschitz stability from such estimates. A new ingredient is the introduction of the cut-off function $\chi_{1}$. See also Lasiecka and Triggiani [19].

Now by using the sharp bounds of Lasiecka, Lions and Triggiani ([18]) in hyperbolic problems under minimal regularity assumptions on their coefficients, we will remove the condition $\mathbf{g}=0$ of Theorem 2.1. In the case of smooth coefficients, such estimates are due to Sakamoto ([24]).

Lemma 2.1. Assume that equations (2.1) do not contain the terms with $\partial_{t} \partial_{k} u_{j}$.

Then there is a constant $C$ depending only on $\Omega, T, P(j), b_{j}, j=1, \ldots, m$ such that any solution $\mathbf{u}$ to (2.1) and (2.2) with the initial conditions $\mathbf{u}=\mathbf{u}_{0}$ and $\partial_{t} \mathbf{u}=\mathbf{u}_{1}$ on $\Omega \times\{0\}$ satisfies the bound

$$
\begin{align*}
& \|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Sigma)  \tag{2.11}\\
& \quad \leq C\left(\left\|\mathbf{u}_{0}\right\|_{(1)}(\Omega)+\left\|\mathbf{u}_{1}\right\|_{2}(\Omega)+\|\mathbf{g}\|_{(1)}(\Sigma)\right)
\end{align*}
$$

when $0<t<T$.
Proof. Let us consider the scalar hyperbolic problem

$$
\left\{\begin{array}{l}
\partial_{t}^{2} v+A v=F \quad \text { in } Q  \tag{2.12}\\
v=v_{0}, \quad \partial_{t} v=v_{1} \quad \text { on } \Omega \times\{0\} \\
v=G \quad \text { on } \Sigma, \quad \text { where } v_{0}=G \text { on } \partial \Omega \times\{0\},
\end{array}\right\}
$$

where $A=-\Sigma \partial_{j}\left(a^{j k} \partial_{k}\right)$ with $a^{j k} \in C^{1}(\bar{Q})$ strictly uniformly positive on $\bar{Q}$. Let $T<T_{0}$. Let $G^{*}$ be an extension of $G$ onto $\Omega \times\left(0, T_{0}\right)$ with $\left\|G^{*}\right\|_{(1)}(\partial \Omega \times$ $\left.\left(0, T_{0}\right)\right) \leq C_{T}\|G\|_{(1)}(\Sigma)$. The existence of $G^{*}$ is guaranteed by extension theorems for Sobolev spaces. Let us extend $F$ onto $\Omega \times\left(T, T_{0}\right)$ as zero. By Theorem 4.1 in [18] for solutions $v^{*}$ to extended problem (2.12), we have

$$
\begin{aligned}
& \left\|v^{*}(\cdot, t)\right\|_{(1)}(\Omega)+\left\|\partial_{t} v^{*}(\cdot, t)\right\|_{2}(\Omega)+\left\|\partial_{\nu} v^{*}\right\|_{2}(\Sigma) \\
& \quad \leq C_{0}\left(\int_{0}^{T}\|F(\cdot, s)\|_{2}(\Omega) d s+\left\|v_{0}\right\|_{(1)}(\Omega)+\left\|v_{1}\right\|_{2}(\Omega)+\left\|G^{*}\right\|_{(1)}\left(\partial \Omega \times\left(0, T_{0}\right)\right)\right)
\end{aligned}
$$

where $C_{0}$ depends only on $\Omega, T_{0}$, the constant of ellipticity of $A$ and on $\left\|a^{l, k}\right\|_{(1)}(Q)$. Hence

$$
\begin{align*}
& \|v(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} v(\cdot, t)\right\|_{2}(\Omega)+\left\|\partial_{\nu} v\right\|_{2}(\Sigma)  \tag{2.13}\\
& \leq C_{0}\left(\int_{0}^{T}\|F(\cdot, s)\|_{2}(\Omega) d s+\left\|v_{0}\right\|_{(1)}(\Omega)+\left\|v_{1}\right\|_{2}(\Omega)+C_{T}\|G\|_{(1)}(\Sigma)\right)
\end{align*}
$$

We will apply bound (2.13) to any of equations (2.1) written in form (2.12) where $A$ is the divergent form of the principal elliptic part of $P(j)$ and $F$ is the sum of $b_{j}$ and of the remainders from the transformation of $P(j)$ into the divergent form. Observe that $\|F(\cdot, s)\|_{2}(\Omega) \leq C\left(\|\mathbf{u}(\cdot, s)\|_{(1)}(\Omega)+\left\|\partial_{t} \mathbf{u}(\cdot, s)\right\|_{2}(\Omega)\right)$ where $C$ depends only on the $C^{1}(\bar{Q})$-norms of the principal coefficients of $P(j)$ and on the $L^{\infty}(Q)$-norms of other coefficients. Summing these bounds over $j=1, \ldots, m$ and applying the triangle inequality, we will have

$$
\begin{aligned}
\|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega) & +\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Sigma) \\
\leq & C_{0}\left(\int_{0}^{T}\left(\|\mathbf{u}(\cdot, s)\|_{(1)}(\Omega)+\left\|\partial_{t} \mathbf{u}(\cdot, s)\right\|_{2}(\Omega)\right) d s\right. \\
& \left.+\left\|\mathbf{u}_{0}\right\|_{(1)}(\Omega)+\left\|\mathbf{u}_{1}\right\|_{2}(\Omega)+C_{T}\|\mathbf{g}\|_{(1)}(\Sigma)\right)
\end{aligned}
$$

for all $t \in(0, T)$. When $T<1 / C_{0}$ by taking supremum of the both parts of the obtained bound for $\mathbf{u}$ and using elementary properties of the integral, one eliminates the integral term. Since smallness of $T$ needed for that is determined only by the coefficients of system (2.1) and since the energy at $t=T$ is bounded by the initial energy and the Dirichlet boundary data, we can repeat this step and in a finite number of steps to exhaust the whole initial interval $(0, T)$. Observe that $T_{0}$ is used to guarantee that the constant $C_{0}$ in (2.13) does not depend on $T$. The proof is complete.

This lemma permits us to eliminate the condition $\mathbf{g}=0$ of Theorem 2.1.
Corollary 2.1. Assume that the operators $P(j)$ do not contain terms with $\partial_{t} \partial_{k} u_{j}$. Let functions $\varphi, \varphi_{1}$ be strongly pseudo-convex and satisfy conditions (1.3), (1.4), and (1.5) as well as other conditions of Theorem 2.1 on $\varphi, \varphi_{1}$.

Then there is a constant $C$ depending only on $\Omega, T, \varphi, P(j), b_{j}, \Gamma$ such that any solution to problem (2.1) and (2.2) satisfies the bound

$$
\begin{align*}
&\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)+\|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Sigma)  \tag{2.14}\\
& \leq C\left(\|\mathbf{g}\|_{(1)}(\Sigma)+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}(\Gamma \times(0, T))\right) .
\end{align*}
$$

Proof. To derive this result from Theorem 2.1 and Lemma 2.1, we introduce a function $\mathbf{u}_{0}^{*} \in H^{1}(\Omega)$ coinciding with $\mathbf{u}$ on $\partial \Omega \times\{0\}$ and such that $\left\|\mathbf{u}_{0}^{*}\right\|_{(1)}(\Omega)$ is bounded by $C\|\mathbf{g}\|_{(1)}(\Sigma)$. Solving the initial value problem for hyperbolic system (2.1) with the initial data $\left(\mathbf{u}^{*}(\cdot, 0), \partial_{t} \mathbf{u}^{*}(\cdot, 0)\right)=\left(\mathbf{u}_{0}^{*}, \mathbf{0}\right)$ and the lateral Dirichlet data $\mathbf{g}$ on $\Sigma$, we obtain a function $\mathbf{u}^{*}$. By Lemma 2.1

$$
\begin{align*}
\left\|\mathbf{u}^{*}(\cdot, t)\right\|_{(1)}(\Omega)+\left\|\partial_{\nu} \mathbf{u}^{*}(\cdot, t)\right\|_{2}(\Omega)+\left\|\partial_{\nu} \mathbf{u}^{*}\right\|_{2}(\Sigma) &  \tag{2.15}\\
& \leq C\|\mathbf{g}\|_{(1)}(\Sigma), \quad 0<t<T .
\end{align*}
$$

Subtracting $\mathbf{u}^{*}$ from $\mathbf{u}$, we obtain for their difference equations (2.1) and the zero lateral Dirichlet data on $\Sigma$. Applying to this difference Theorem 2.1 and combining its bound with (2.15) by the triangle inequality, we complete the proof.

Now we discuss cases (1) and (2) in the $n$-dimensional case:
(1) $\Omega \subset B(0, R)$ and $\Gamma \subset \partial \Omega$,
(2) $\Omega \subset\left\{-h<x_{n}<0, x_{1}^{2}+\cdots+x_{n-1}^{2}<r^{2}\right\}$,
which are already introduced before Theorem 1.2 (the three dimensional case).
Corollary 2.2. Assume that the coefficients $a=a(j), j=1, \ldots, m$ of the principal parts $P(j)=\square_{a(j)}$ of system (2.1) satisfy conditions (1.7).

If in case (1) $R<\theta T / 2$ and $\Gamma=\partial \Omega$ and in case (2) $\Gamma$ contains $\{0 \leq$ $\left.x \cdot \nu-\beta \nu_{n}\right\} \cap \partial \Omega, r^{2}+h^{2}+2 h \beta<\theta^{2} T^{2} / 4$, then bound (2.14) holds.

Proof. We will apply Corollary 2.1 with

$$
\varphi=e^{\sigma \psi}-1
$$

and

$$
\psi(x, t)=x_{1}^{2}+\cdots+x_{n-1}^{2}+\left(x_{n}-\beta\right)^{2}-\theta^{2}(t-T / 2)^{2}-s .
$$

It is known that under conditions (1.7) with respect to the principal coefficients the function $\varphi=\exp (\sigma \psi)$ for large $\sigma$ is strongly pseudo-convex in $\bar{Q} \backslash\{(\beta, T / 2)\}$ with respect to the wave operators $\square_{a(j)}$ ([14, Section 3.4]).

In case (1) we let $\beta=0$ and will make use also of the funcion $\varphi_{1}=e^{\lambda \psi_{1}}$, $\psi_{1}(x, t)=\psi\left(x_{1}, \ldots, x_{n-1}, x_{n}-\beta_{0}, t\right)$ with small $\beta_{0}$. Then condition (1.3) is satisfied with any $s>R^{2}-\theta^{2} T^{2} / 4$ and condition (1.4) is satisfied with any $s<0$. Since $R<\theta T / 2$ we can satisfy the both conditions (choosing for example $\left.s=(1 / 2)\left(R^{2}-\theta^{2} T^{2} / 4\right)\right)$.

In case (2) we let $s=\beta^{2}-\delta, \delta>0$. Then the function $\varphi$ is strongly pseudo-convex on the whole $\bar{Q}$ and the function $\varphi_{1}$ is not needed. But to
comply formally with Theorem 2.1 we can choose as $x^{0}, x^{1}$ any points of $\Omega$ such that $\varphi\left(x^{1}, T / 2\right)<\varphi\left(x^{0}, T / 2\right)$ and let $\varphi_{1}=\varphi$. Condition (1.3) follows from the inequality $x_{1}^{2}+\cdots+x_{n-1}^{2}+\left(x_{n}-\beta\right)^{2}-\theta^{2} T^{2} / 4-\beta^{2}+\delta<0$ for all $x \in \bar{\Omega}$ which is a corollary of the inequality $r^{2}+h^{2}+2 h \beta<\theta^{2} T^{2} / 4$ when $\delta$ is small. Condition (1.4) is satisfied because $0<\delta=\beta^{2}-\left(\beta^{2}-\delta\right) \leq \psi(x, T / 2)$.

So in the both cases Corollary 2.2 follows from Corollary 2.1.
It is not difficult to see that conditions (1.7) and $r^{2}+h^{2}+2 h \beta<\theta^{2} T^{2} / 4$ are compatible and mean that $T$ is large (and $\theta$ is small, but it suffices to let $\theta=T^{-3 / 4}$ to satisfy the both conditions). Condition (1.7) can be satisfied when $0<a_{j}+(1 / 2) x \cdot \nabla a_{j}-(\beta / 2) \partial_{n} a_{j}$ on $\bar{Q}$.

## 3. Proofs of estimates for the elasticity system

To prove Theorems 1.1 and 1.2, we will extend system (1.1) for three unknown functions $u_{1}, u_{2}, u_{3}$ to a new one for four unknown functions by introducing $v=\nabla^{T} \mathbf{u}$. We refer to Eller, Isakov, Nakamura and Tataru [7], Ikehata, Nakamura and Yamamoto [12].

Lemma 3.1. Let $v=\nabla^{T} \mathbf{u}$. If $\mathbf{u}$ solves (1.1), then

$$
\left\{\begin{array}{c}
\rho \partial_{t}^{2} v-(\lambda+2 \mu) \Delta v+A_{1 ; 1}(v, \mathbf{u})=0  \tag{3.1}\\
\rho \partial_{t}^{2} \mathbf{u}-\mu \Delta \mathbf{u}+A_{1 ; 2}(v, \mathbf{u})=0 \quad \text { in } Q
\end{array}\right\}
$$

where $A_{1 ; 1}$ and $A_{1 ; 2}$ are (matrix) linear partial differential operators with measurable and bounded coefficients in $Q$.

This lemma and the results of Section 2 imply Theorems 1.1 and 1.2 for Cauchy problem (1.1) and (1.2).

Next we have to prove Theorems 1.3 and 1.4. We recall that we take

$$
\begin{equation*}
\mathbf{g}=0 \tag{3.2}
\end{equation*}
$$

For simplicity we set

$$
\begin{align*}
\mathcal{L} \mathbf{u}= & \mu\left(\Delta \mathbf{u}+\nabla\left(\nabla^{T} \mathbf{u}\right)\right)-\nabla\left(\lambda \nabla^{T} \mathbf{u}\right) \\
& -\sum_{j=1}^{3} \nabla \mu \cdot\left(\nabla u_{j}+\partial_{j} \mathbf{u}\right) \mathbf{e}_{j} \tag{3.3}
\end{align*}
$$

and $\mathbf{a}^{T}$ denotes the transpose of a vector a under consideration.
For the proof of Theorems 1.3 and 1.4, it is sufficient to prove
Lemma 3.2. On $\partial \Omega$, we represent $\nabla \mathbf{u}$ and $\partial_{\nu}\left(\nabla^{T} \mathbf{u}\right)$ by linear functions of $\sigma(\mathbf{u}) \nu$ and $\nabla(\sigma(\mathbf{u}) \nu)$ with bounded coefficients.

Proof. The proof is modification of an argument in [12] where $\sigma(\mathbf{u}) \nu=0$ is assumed. Here we give the proof for completeness.

First Step. We set

$$
\nabla \mathbf{u}=\left(\frac{\partial u_{i}}{\partial x_{j}}\right)_{1 \leq i, j \leq 3}=\left(\begin{array}{c}
\left(\nabla u_{1}\right)^{T}  \tag{3.4}\\
\left(\nabla u_{2}\right)^{T} \\
\left(\nabla u_{3}\right)^{T}
\end{array}\right) .
$$

Then, under assumption (3.2), we have

$$
\begin{equation*}
\nabla \mathbf{u}=\{(\nabla \mathbf{u}) \nu\} \nu^{T} \quad \text { on } \quad \partial \Omega \times(0, T) \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\nabla^{T} \mathbf{u}=(\nabla \mathbf{u}) \nu \cdot \nu \quad \text { on } \quad \partial \Omega \times(0, T) \tag{3.6}
\end{equation*}
$$

In fact, setting $\mathbf{u}=\left(u_{1}, u_{2}, u_{3}\right)^{T}$ and $\nu=\left(\nu_{1}, \nu_{2}, \nu_{3}\right)^{T}$, we see that condition (3.2) implies

$$
\nabla u_{i}=\left(\begin{array}{c}
\frac{\partial u_{i}}{\partial x_{1}} \\
\frac{\partial u_{i}}{\partial x_{2}} \\
\frac{\partial u_{i}}{\partial x_{3}}
\end{array}\right)=\left(\nabla u_{i} \cdot \nu\right) \nu=\left(\frac{\partial u_{i}}{\partial \nu}\right) \nu, \quad 1 \leq i \leq 3 .
$$

Therefore we have

$$
\nabla \mathbf{u}=\left(\begin{array}{l}
\left(\nabla u_{1}\right)^{T}  \tag{3.7}\\
\left(\nabla u_{2}\right)^{T} \\
\left(\nabla u_{3}\right)^{T}
\end{array}\right)=\left(\begin{array}{l}
\left(\nabla u_{1} \cdot \nu\right) \nu^{T} \\
\left(\nabla u_{2} \cdot \nu\right) \nu^{T} \\
\left(\nabla u_{3} \cdot \nu\right) \nu^{T}
\end{array}\right) .
$$

By definition (3.4), this means (3.5). Moreover by $\nu \nu^{T}=1$, we have

$$
(\nabla \mathbf{u}) \nu=\left(\begin{array}{c}
\left(\nabla u_{1} \cdot \nu\right) \nu^{T} \\
\left(\nabla u_{2} \cdot \nu\right) \nu^{T} \\
\left(\nabla u_{3} \cdot \nu\right) \nu^{T}
\end{array}\right) \nu=\left(\begin{array}{c}
\left(\nabla u_{1} \cdot \nu\right) \\
\left(\nabla u_{2} \cdot \nu\right) \\
\left(\nabla u_{3} \cdot \nu\right)
\end{array}\right),
$$

and so

$$
(\nabla \mathbf{u}) \nu \cdot \nu=\left(\nabla u_{1} \cdot \nu\right) \nu_{1}+\left(\nabla u_{2} \cdot \nu\right) \nu_{2}+\left(\nabla u_{3} \cdot \nu\right) \nu_{3}=\operatorname{Trace} \nabla \mathbf{u}=\nabla^{T} \mathbf{u}
$$

by (3.7). Thus we see (3.5) and (3.6).
Second Step. In this step, we will prove that

$$
\begin{equation*}
\nabla \mathbf{u}=S(\sigma(\mathbf{u}) \nu) \quad \text { on } \quad \partial \Omega \times(0, T) \tag{3.8}
\end{equation*}
$$

for some $3 \times 3$ matrix $S \in C(\partial \Omega \times[0, T])^{3 \times 3}$.
Proof of (3.8). We define a $3 \times 3$ matrix $B=B(x, t)$ by

$$
\begin{equation*}
B(x, t) \mathbf{a}=\lambda(x, t)(\mathbf{a} \cdot \nu(x)) \nu(x)+2 \mu(x, t)\left\{\operatorname{Sym}\left(\mathbf{a} \nu(x)^{T}\right)\right\} \nu(x) \tag{3.9}
\end{equation*}
$$

for $\mathbf{a} \in \mathbb{R}^{3}$. Here and henceforth, for square matrices $A=\left(a_{i j}\right)_{1 \leq i, j \leq 3}$ and $B=\left(b_{i j}\right)_{1 \leq i, j \leq 3}$, we set $\operatorname{Sym} A=(1 / 2)\left(A+A^{T}\right)$ and $A \cdot B=\sum_{i, j=1}^{3} a_{i j} b_{i j}$, $|A|^{2}=\sum_{i, j=1}^{3} a_{i j}^{2}$. Then

$$
\begin{equation*}
B=B(x, t) \text { is invertible for all }(x, t) \in \partial \Omega \times(0, T) \tag{3.10}
\end{equation*}
$$

In fact, since $C \nu \cdot \mathbf{a}=C \cdot\left(\mathbf{a} \nu^{T}\right)$ and $(1 / 2)\left(C+C^{T}\right) \cdot C=\left|(1 / 2)\left(C+C^{T}\right)\right|^{2}$ for an $n \times n$ matrix $C$, we have

$$
\left(\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right) \nu \cdot \mathbf{a}=\left(\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right) \cdot\left(\mathbf{a} \nu^{T}\right)=\left|\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right|^{2} .
$$

Hence

$$
\begin{align*}
(B \mathbf{a} \cdot \mathbf{a}) & =\lambda|\mathbf{a} \cdot \nu|^{2}+2 \mu\left|\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right|^{2} \\
& =\lambda|\operatorname{Trace} A|^{2}+2 \mu|A|^{2} . \tag{3.11}
\end{align*}
$$

Here we set

$$
A=\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)
$$

and

$$
\begin{equation*}
C=A-\frac{\operatorname{Trace} A}{3} I_{3}, \tag{3.12}
\end{equation*}
$$

where $I_{3}$ is the $3 \times 3$ identity matrix. Then Trace $C=0$, so that

$$
\begin{equation*}
C \cdot I_{3}=0 \tag{3.13}
\end{equation*}
$$

by the identity $C \cdot I_{3}=$ Trace $C$. Therefore (3.11) through (3.13) imply

$$
\begin{aligned}
B \mathbf{a} \cdot \mathbf{a} & =\lambda|\operatorname{Trace} A|^{2}+2 \mu\left|\frac{\operatorname{Trace} A}{3} I_{3}+C\right|^{2} \\
& =\lambda|\operatorname{Trace} A|^{2}+2 \mu\left(\left|\frac{\operatorname{Trace} A}{3} I_{3}\right|^{2}+|C|^{2}+2 \frac{\operatorname{Trace} A}{3} I_{3} \cdot C\right) \\
& =\frac{3 \lambda+2 \mu}{3}|\operatorname{Trace} A|^{2}+2 \mu|C|^{2} \geq \frac{\delta_{0}}{3}|\operatorname{Trace} A|^{2}+\delta_{0}|C|^{2},
\end{aligned}
$$

where $\delta_{0}>0$ is a constant. By (3.12), we have $A=C+($ Trace $A / 3) I_{3}$, so that $\delta_{0}|A|^{2}=\left(\delta_{0} / 3\right) \mid$ Trace $\left.A\right|^{2}+\delta_{0}|C|^{2}$ by (3.13). Therefore

$$
B \mathbf{a} \cdot \mathbf{a} \geq \delta_{0}\left|\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right|^{2} .
$$

Moreover we have

$$
\begin{aligned}
\left|\operatorname{Sym}\left(\mathbf{a} \nu^{T}\right)\right|^{2} & =\frac{1}{4}\left(\left|\mathbf{a} \nu^{T}\right|^{2}+2\left(\mathbf{a} \nu^{T}\right) \cdot\left(\nu \mathbf{a}^{T}\right)+\left|\nu \mathbf{a}^{T}\right|^{2}\right) \\
& =\frac{1}{4}\left(|\mathbf{a}|^{2}+2|\mathbf{a} \cdot \nu|^{2}+|\mathbf{a}|^{2}\right) \geq \frac{1}{2}|\mathbf{a}|^{2}
\end{aligned}
$$

so that

$$
\begin{equation*}
B \mathbf{a} \cdot \mathbf{a} \geq \frac{\delta_{0}}{2}|\mathbf{a}|^{2} \quad \text { on } \quad \partial \Omega \times(0, T) \tag{3.14}
\end{equation*}
$$

Moreover direct calculations verify that $B \mathbf{a} \cdot \mathbf{b}=B \mathbf{b} \cdot \mathbf{a}$ for every $\mathbf{a}, \mathbf{b} \in \mathbb{R}^{3}$, which means that $B$ is a symmetric matrix. Therefore (3.14) implies (3.10).

On the other hand, by (3.5) and (3.6) we have

$$
\begin{aligned}
B((\nabla \mathbf{u}) \nu) & =\lambda\{(\nabla \mathbf{u}) \nu \cdot \nu\} \nu+2 \mu\left\{\operatorname{Sym}\left(((\nabla \mathbf{u}) \nu) \nu^{T}\right)\right\} \nu \\
& =\lambda\left(\nabla^{T} \mathbf{u}\right) \nu+2 \mu\{\operatorname{Sym}(\nabla \mathbf{u})\} \nu=\sigma(\mathbf{u}) \nu \quad \text { on } \quad \partial \Omega \times(0, T) .
\end{aligned}
$$

In view of (3.5) and (3.10), the proof of (3.8) is complete.
Third Step. Since $\mathbf{u} \in C^{2}(\bar{Q})^{n}$, we see from (3.2) and (3.8) that $\mathcal{L} \mathbf{u}=$ $\mu\left(\Delta \mathbf{u}+\nabla\left(\nabla^{T} \mathbf{u}\right)\right)-\nabla\left(\lambda \nabla^{T} \mathbf{u}\right)+F_{1}(S(\sigma(\mathbf{u}) \nu))$ on $\partial \Omega \times(0, T)$. Here and henceforth $F_{i}$ denote linear maps with continuous coefficients. Therefore we obtain

$$
\begin{equation*}
\mu \Delta \mathbf{u}+(\lambda+\mu) \nabla\left(\nabla^{T} \mathbf{u}\right)=F_{2}(\sigma(\mathbf{u}) \nu) \quad \text { on } \quad \partial \Omega \times(0, T) . \tag{3.15}
\end{equation*}
$$

Since $\partial \Omega$ is of $C^{2}$-class, for any $x^{0}=\left(x_{1}^{0}, x_{2}^{0}, x_{3}^{0}\right) \in \partial \Omega$, we can take neighbourhoods $\mathcal{V}$ in $\mathbb{R}^{3}$ of $x^{0}$ and $\mathcal{U}$ in $\mathbb{R}^{2}$ of $\left(x_{1}^{0}, x_{2}^{0}\right)$, a function $\chi=\chi\left(x_{1}, x_{2}\right) \in C^{2}(\mathcal{U})$ such that

$$
\begin{equation*}
\left(x_{1}, x_{2}, x_{3}\right) \in \mathcal{V} \cap \partial \Omega \quad \text { if and only if } \quad x_{3}=\chi\left(x_{1}, x_{2}\right) \tag{3.16}
\end{equation*}
$$

We introduce a new coordinate $\eta=\eta(x)=\left(\eta_{1}, \eta_{2}, \eta_{3}\right)$ by

$$
\begin{equation*}
\eta_{1}=x_{1}, \quad \eta_{2}=x_{2}, \quad \eta_{3}=x_{3}-\chi\left(x_{1}, x_{2}\right) \tag{3.17}
\end{equation*}
$$

for $\left(x_{1}, x_{2}\right) \in \mathcal{U}$. We define a set $\mathcal{W}$ of $\left(\eta_{1}, \eta_{2}, \eta_{3}\right)$ by $\mathcal{W}=\eta(\mathcal{V} \cap \partial \Omega) \subset$ $\left\{\left(\eta_{1}, \eta_{2}, 0\right) ; \eta_{1}, \eta_{2} \in \mathbb{R}\right\}$. Henceforth we locally regard $\mathbf{u}=\mathbf{u}\left(x_{1}, x_{2}, x_{3}\right)$ as a function in $\left(\eta_{1}, \eta_{2}, \eta_{3}\right) \in \eta(\mathcal{V})$. Then boundary conditions (3.2) and (3.8) imply

$$
\begin{align*}
& \mathbf{u}=0, \quad \frac{\partial \mathbf{u}}{\partial \eta_{j}}=F_{3}(\sigma(\mathbf{u}) \nu)  \tag{3.18}\\
& \frac{\partial^{2} \mathbf{u}}{\partial \eta_{i} \partial \eta_{j}}=F_{4}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)), \quad i=1,2, \quad j=1,2,3, \quad \text { in } \mathcal{W} .
\end{align*}
$$

Then noting that $\partial_{3} \chi=-1$,

$$
\frac{\partial \mathbf{u}}{\partial x_{i}}=\frac{\partial \mathbf{u}}{\partial \eta_{i}}-\partial_{i} \chi \frac{\partial \mathbf{u}}{\partial \eta_{3}}, \quad i=1,2, \quad \frac{\partial \mathbf{u}}{\partial x_{3}}=\frac{\partial \mathbf{u}}{\partial \eta_{3}}
$$

we see

$$
\begin{aligned}
\frac{\partial^{2} \mathbf{u}}{\partial x_{i} \partial x_{j}}= & \frac{\partial^{2} \mathbf{u}}{\partial \eta_{i} \partial \eta_{j}}-\partial_{j} \chi \frac{\partial^{2} \mathbf{u}}{\partial \eta_{i} \partial \eta_{3}}-\left(\partial_{j} \partial_{i} \chi\right) \frac{\partial \mathbf{u}}{\partial \eta_{3}} \\
& -\partial_{i} \chi \frac{\partial^{2} \mathbf{u}}{\partial \eta_{j} \partial \eta_{3}}+\left(\partial_{i} \chi\right)\left(\partial_{j} \chi\right) \frac{\partial^{2} \mathbf{u}}{\partial \eta_{3}^{2}}, \quad 1 \leq i, j \leq 2 \\
\frac{\partial^{2} \mathbf{u}}{\partial x_{3}^{2}}= & \frac{\partial^{2} \mathbf{u}}{\partial \eta_{3}^{2}}, \quad \frac{\partial^{2} \mathbf{u}}{\partial x_{i} \partial x_{3}}=\frac{\partial^{2} \mathbf{u}}{\partial \eta_{i} \partial \eta_{3}}-\partial_{i} \chi \frac{\partial^{2} \mathbf{u}}{\partial \eta_{3}^{2}}, \quad i=1,2
\end{aligned}
$$

Therefore (3.18) implies

$$
\begin{equation*}
\frac{\partial^{2} \mathbf{u}}{\partial x_{i} \partial x_{j}}=\partial_{i} \chi \partial_{j} \chi \frac{\partial^{2} \mathbf{u}}{\partial \eta_{3}^{2}}+F_{5}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)), \quad 1 \leq i, j \leq 3 \tag{3.19}
\end{equation*}
$$

in $\mathcal{W}$. We substitute (3.19) into (3.15), and we obtain

$$
\begin{array}{r}
\mu|\nabla \chi|^{2} \frac{\partial^{2} u_{i}}{\partial \eta_{3}^{2}}+(\lambda+\mu)\left(\partial_{i} \chi \partial_{1} \chi \frac{\partial^{2} u_{1}}{\partial \eta_{3}^{2}}+\partial_{i} \chi \partial_{2} \chi \frac{\partial^{2} u_{2}}{\partial \eta_{3}^{2}}+\partial_{i} \chi \partial_{3} \chi \frac{\partial^{2} u_{3}}{\partial \eta_{3}^{2}}\right) \\
=F_{6}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)), \quad 1 \leq i \leq 3
\end{array}
$$

in $\mathcal{W}$. We can rewrite the above equalities as

$$
\begin{equation*}
D(\eta, t)\left(\frac{\partial^{2} u_{1}}{\partial \eta_{3}^{2}}, \frac{\partial^{2} u_{2}}{\partial \eta_{3}^{2}}, \frac{\partial^{2} u_{3}}{\partial \eta_{3}^{2}}\right)^{T}=F_{7}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)) \quad \text { in } \quad \mathcal{W} \times(0, T) \tag{3.20}
\end{equation*}
$$

where we define a $3 \times 3$ matrix $D=D(\eta, t)$ by

$$
D=D(\eta, t)=(\lambda+\mu)(\nabla \chi)(\nabla \chi)^{T}+\mu|\nabla \chi|^{2} I_{3} .
$$

For any $\mathbf{a} \in \mathbb{R}^{3}$, we have

$$
\begin{align*}
D \mathbf{a} \cdot \mathbf{a} & =(\lambda+\mu)\left((\nabla \chi)(\nabla \chi)^{T} \mathbf{a} \cdot \mathbf{a}\right)+\mu|\nabla \chi|^{2}(\mathbf{a} \cdot \mathbf{a}) \\
& =(\lambda+\mu)(\mathbf{a} \cdot \nabla \chi)^{2}+\mu|\nabla \chi|^{2}|\mathbf{a}|^{2} . \tag{3.21}
\end{align*}
$$

We see from $\mu>0$ and $3 \lambda+2 \mu>0$ on $\bar{Q}$ that

$$
\begin{equation*}
\lambda(x, t)+2 \mu(x, t)>\delta_{0}, \quad(x, t) \in \partial \Omega \times(0, T) \tag{3.22}
\end{equation*}
$$

where $\delta_{0}>0$ is a constant.
Since $\mu \in C(\bar{Q})$, by (3.22) we can choose a sufficiently small $\varepsilon>0$ such that

$$
\begin{equation*}
\lambda(x, t)+(2-\varepsilon) \mu(x, t)>0, \quad \mu(x, t) \geq \frac{\delta_{0}}{2}, \quad(x, t) \in \partial \Omega \times(0, T) . \tag{3.23}
\end{equation*}
$$

Applying Schwarz's inequality in (3.21), by (3.23), we obtain

$$
\begin{align*}
D \mathbf{a} \cdot \mathbf{a}= & (\lambda+\mu)(\mathbf{a} \cdot \nabla \chi)^{2}+(\mu-\mu \varepsilon)|\nabla \chi|^{2}|\mathbf{a}|^{2}+\mu \varepsilon|\nabla \chi|^{2}|\mathbf{a}|^{2} \\
= & (\lambda+(2-\varepsilon) \mu)(\mathbf{a} \cdot \nabla \chi)^{2} \\
& +\mu(1-\varepsilon)\left(|\nabla \chi|^{2}|\mathbf{a}|^{2}-(\nabla \chi \cdot \mathbf{a})^{2}\right)+\mu \varepsilon|\nabla \chi|^{2}|\mathbf{a}|^{2}  \tag{3.24}\\
\geq & \mu \varepsilon|\nabla \chi|^{2}|\mathbf{a}|^{2} \geq \frac{\delta_{0} \varepsilon}{2}|\nabla \chi|^{2}|\mathbf{a}|^{2} \geq \frac{\delta_{0} \varepsilon}{2}|\mathbf{a}|^{2} .
\end{align*}
$$

At the last inequality, we use $|\nabla \chi|^{2}=1+\left(\partial_{1} \chi\right)^{2}+\left(\partial_{2} \chi\right)^{2} \geq 1$. By the definition, $D(\eta, t)$ is symmetric, inequality (3.24) implies that $D=D(\eta, t)$ is invertible in $\mathcal{W} \times(0, T)$. Therefore (3.20) yields

$$
\frac{\partial^{2} \mathbf{u}}{\partial \eta_{3}^{2}}=F_{8}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)) \quad \text { in } \quad \mathcal{W} \times(0, T)
$$

We combine (3.19) to obtain

$$
\frac{\partial^{2} \mathbf{u}}{\partial x_{i} \partial x_{j}}(x, t)=F_{9}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)), \quad x \in \mathcal{V} \cap \partial \Omega, 0<t<T, \leq i, j \leq 3
$$

Therefore, since $x^{0} \in \mathcal{V}$ is an arbitrary point of $\partial \Omega$ and $\partial \Omega$ is compact,

$$
\partial_{\nu}\left(\nabla^{T} \mathbf{u}\right)=F_{9}(\sigma(\mathbf{u}) \nu, \nabla(\sigma(\mathbf{u}) \nu)) \quad \text { on } \quad \partial \Omega \times(0, T),
$$

where $F_{9}$ is a linear map whose coefficients are bounded on $\partial \Omega \times(0, T)$. Thus in view of Theorem 1.1, the proof of Theorem 1.3 is complete.
4. Lipschitz stability for a principally diagonal hyperbolic system on the basis of a directly derived Carleman estimate

Theorems 1.1 through 1.4 are proved on the basis of Carleman estimates with boundary data by Tataru [27], where Carleman estimates are shown in a general setting. On the other hand, we can establish a similar Carleman estimate more directly for an operator $\square_{a} \equiv a \partial_{t}^{2}-\Delta$. This direct derivation is shown in Lavrent'ev, Romanov and Shishat•skiĭ [20] for the operator $\square$ (i.e. $a \equiv 1$ ). In this section, we modify the argument in $[20]$ for $\square_{a}$. This direct way is less transparent than [27], while it admits less regular boundary $\partial \Omega$. On the other hand, the direct way gives a worse condition for $T$ than in Sections 1 and 2 , although in the case of constant $a$, the results in this section coincide with the results by Tataru's Carleman estimates and by the multiplier method.

In this section, we treat general spatial dimensions again except for Theorem 4.3. First for the statement of the Carleman estimate directly derived, we take $x_{0} \in \mathbb{R}^{n}$ with $x_{0} \notin \bar{\Omega}$ and let us assume

$$
\begin{equation*}
a(x)>0, \quad 2+\frac{\left(\nabla a(x) \cdot\left(x-x_{0}\right)\right)}{a(x)}>0, \quad x \in \bar{\Omega} . \tag{4.1}
\end{equation*}
$$

We take $\beta>0$ such that

$$
\left\{\begin{array}{l}
\beta<\frac{4}{T|\nabla a(x)|},  \tag{4.2}\\
\beta<\frac{4 a(x)+2\left(\nabla a(x) \cdot\left(x-x_{0}\right)\right)}{T|\nabla a(x)|(1+a(x))+4 a(x)^{2}}, \\
\beta<\frac{1}{a(x)} \quad \text { on } \bar{Q} .
\end{array}\right\} .
$$

Remark 4.1. In the case of $a=$ constant, condition (4.2) is reduced to $\beta<(1 / a)$.

Moreover we set

$$
\begin{equation*}
\varphi(x, t)=\left|x-x_{0}\right|^{2}-\beta\left(t-\frac{T}{2}\right)^{2} \tag{4.3}
\end{equation*}
$$

and

$$
\begin{equation*}
Q(\delta)=\left\{(x, t): x \in \Omega, \varphi(x, t) \geq \delta^{2}\right\} \tag{4.4}
\end{equation*}
$$

with $\delta>0$,

$$
\begin{equation*}
\Gamma_{+}=\left\{x \in \partial \Omega:\left(\left(x-x_{0}\right) \cdot \nu(x)\right)>0\right\}, \quad \Gamma_{-}=\partial \Omega \backslash \Gamma_{+} \tag{4.5}
\end{equation*}
$$

Now we can state a Carleman estimate.
Theorem 4.1. For $\delta>0$, there exists a constant $C=C(\delta)>0$ such that

$$
\begin{align*}
& \int_{Q(\delta)}\left|\square_{a} u\right|^{2} e^{2 \tau \varphi} d x d t+\int_{\left(\Gamma_{+} \times(0, T)\right) \cap \partial Q(\delta)} \tau e^{2 \tau \varphi}\left|\partial_{\nu} u\right|^{2} d \sigma  \tag{4.6}\\
& \quad \geq C \tau \int_{Q(\delta)}\left(|\nabla u|^{2}+\left|\partial_{t} u\right|^{2}\right) e^{2 \tau \varphi} d x d t+C \tau^{3} \int_{Q(\delta)} u^{2} e^{2 \tau \varphi} d x d t
\end{align*}
$$

for large $\tau>0$ provided that

$$
\left\{\begin{array}{lll}
u \in H_{0}^{1}(Q(\delta)) \cap H^{2}(Q(\delta)), & \nabla u=0 \quad \text { on } \quad \partial Q(\delta) \backslash(\partial \Omega \times(0, T))  \tag{4.7}\\
\partial_{t} u=0 \quad \text { on } \quad \partial Q(\delta) .
\end{array}\right.
$$

We consider

$$
\begin{equation*}
a_{j}(x) \partial_{t}^{2} u_{j}-\Delta u_{j}+b_{j}\left(x, t ; \nabla_{x, t} \mathbf{u}, \mathbf{u}\right)=0, \quad \text { in } \quad Q, 1 \leq j \leq m \tag{4.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbf{u}=\mathbf{g} \quad \text { on } \quad \Gamma . \tag{4.9}
\end{equation*}
$$

Here $\mathbf{u}=\left(u_{1}, \ldots, u_{m}\right)$ and $a_{j}, 1 \leq j \leq m$ satisfy (4.1), $b_{j}$ are linear functions of $\nabla \mathbf{u}$ and $\mathbf{u}$ whose coefficients are in $L^{\infty}(Q)$ if the order of differentiation is 1 and the coefficient of the zeroth order term is in $L^{n+1}(Q)$.

On the basis of Theorem 4.1, we can prove
Theorem 4.2. Let $\Omega$ be a bounded domain in $\mathbb{R}^{n}$ whose boundary $\partial \Omega$ is of piecewise $C^{1}$. We assume

$$
\begin{equation*}
T>2 \max _{x \in \bar{\Omega}} \frac{\left|x-x_{0}\right|}{\sqrt{\beta}} . \tag{4.10}
\end{equation*}
$$

Here $\beta>0$ satisfies (4.2) for all $a_{1}, \ldots, a_{m}$. Then there exists a constant $C>0$ depending on $\Omega, T, x_{0}, a_{j}, b_{j}, 1 \leq j \leq m$, such that any solution $\mathbf{u}$ to (4.8) and (4.9) satisfies the bound

$$
\begin{aligned}
& \|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega) \\
& \quad \leq C\left(\|\mathbf{g}\|_{(1)}(\Sigma)+\left\|\partial_{\nu} \mathbf{u}\right\|_{2}\left(\Gamma_{+} \times(0, T)\right)\right), \quad 0 \leq t \leq T .
\end{aligned}
$$

Remark 4.2. We set

$$
\gamma(a)(x, t)=\min \left\{\frac{4}{T|\nabla a(x)|}, \frac{4 a(x)+2\left(\nabla a(x) \cdot\left(x-x_{0}\right)\right)}{T|\nabla a|(1+a)+4 a^{2}}, \frac{1}{a}\right\} .
$$

Then (4.2) is rewritten as

$$
\begin{equation*}
\beta<\min _{(x, t) \in \bar{Q}} \gamma(a)(x, t) . \tag{4.2}
\end{equation*}
$$

In terms of $\gamma(a)$, we can rewrite (4.10) by

$$
T>2 \max _{(x, t) \in \bar{Q}, 1 \leq j \leq m} \frac{\left|x-x_{0}\right|}{\sqrt{\gamma\left(a_{j}\right)(x, t)}}
$$

In the case where $a_{j}, 1 \leq j \leq m$ are constants, then condition (4.10) for $T>0$ follows from

$$
\begin{equation*}
T>2 \max _{x \in \bar{\Omega}, 1 \leq j \leq m} \sqrt{a_{j}}\left|x-x_{0}\right| \tag{4.10}
\end{equation*}
$$

which is same as required by the multiplier method (e.g. Isakov [14], Komornik [17], Lop Fat Ho [22], Powell [23]). Therefore our result can generalize observability with constant coefficients by the multiplier method. We notice that the multiplier method can not treat $b_{j}$ with terms of derivatives of the first order.

By the same argument as in Section 3, Theorem 4.2 yields
Theorem 4.3. Let $\Omega \subset \mathbb{R}^{3}$ be a bounded domain and let its boundary $\partial \Omega$ be of piecewise $C^{2}$. We set

$$
\begin{equation*}
a_{1}(x)=\frac{\rho(x)}{\mu(x)}, \quad a_{2}(x)=\frac{\rho(x)}{2 \mu(x)+\lambda(x)}, \quad x \in \Omega . \tag{4.11}
\end{equation*}
$$

In (1.1) we assume that $\rho, \lambda, \mu \in C^{3}(\bar{Q})$ are independent of $t$, and $\mathbf{g}=0$ in (1.2). For $x_{0} \notin \bar{\Omega}$, we assume that $a_{1}$ and $a_{2}$ satisfy (4.1),

$$
\begin{equation*}
T>2 \max _{x \in \bar{\Omega}} \frac{\left|x-x_{0}\right|}{\sqrt{\beta}} \tag{4.12}
\end{equation*}
$$

where $\beta>0$ satisfies (4.2) for $a_{1}$ and $a_{2}$. Then there exists a constant $C>0$ depending on $\Omega, T, x_{0}, \rho, \mu, \lambda$, such that any solution $\mathbf{u}$ to (1.1) and (1.2)' satisfies the bound

$$
\begin{align*}
& \|\mathbf{u}(\cdot, t)\|_{(1)}(\Omega)+\left\|\partial_{t} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega) \\
& \quad+\left\|\left(\nabla^{T} \mathbf{u}\right)(\cdot, t)\right\|_{(1)}(\Omega)+\left\|\partial_{t} \nabla^{T} \mathbf{u}(\cdot, t)\right\|_{2}(\Omega)  \tag{4.13}\\
& \quad \leq C\left\{\|\mathbf{h}\|_{2}\left(\Gamma_{+} \times(0, T)\right)+\|\nabla \mathbf{h}\|_{2}\left(\Gamma_{+} \times(0, T)\right)\right\}, \quad 0 \leq t \leq T .
\end{align*}
$$

A similar bound holds for solutions to (1.1) and (1.2).

The rest part of this section is devoted to the proofs of Theorems 4.1 and 4.2.

Derivation of Theorem 4.2 from Theorem 4.1. The deduction is similar to the proof of Theorem 2.1. By $x_{0} \notin \bar{\Omega}$, we note that we can take $\varphi_{1}=\varphi$. In terms of Lemma 2.1, we may assume that $\mathbf{g}=0$. We set $t_{0}=T / 2$ in (4.3), that is,

$$
\begin{equation*}
\varphi(x, t)=\left|x-x_{0}\right|^{2}-\beta\left(t-\frac{T}{2}\right)^{2} \tag{4.14}
\end{equation*}
$$

where $\beta>0$ is chosen later. For the proof, it is sufficient to verify that $\varphi$ defined by (4.14) satisfies (1.3) and (1.4) with $\beta>0$ satisfying (4.2). In fact, condition (1.4) is straightforward from $x_{0} \notin \bar{\Omega}$ and (1.3) follows directly from (4.10).

Proof of Theorem 4.1. We will extend the proof in [20] (pp. 123-128) where $a(x) \equiv 1$ is considered. For convenience, we denote:

$$
\begin{align*}
& x=\left(x_{1}, \ldots, x_{n}\right), \quad x_{n+1}=t, \quad A_{1}=\cdots=A_{n}=1, \quad A_{n+1}=-a, \\
& w(x, t)=e^{\tau \varphi(x, t)} \tag{4.15}
\end{align*}
$$

and will make use of the function

$$
\begin{equation*}
\psi(x, t)=4-2 \beta|\nabla a(x)|\left|t-\frac{T}{2}\right|-2 n-2 \beta a(x)-\varepsilon . \tag{4.16}
\end{equation*}
$$

Letting

$$
\begin{equation*}
v=w u \tag{4.17}
\end{equation*}
$$

we obtain

$$
\begin{align*}
w^{2}\left(\square_{a} u\right)^{2}= & \left\{\sum_{j=1}^{n+1} A_{j}\left\{\partial_{j}^{2} v-2 \tau \partial_{j} \varphi \partial_{j} v+\left(\tau^{2}\left(\partial_{j} \varphi\right)^{2}-\tau \partial_{j}^{2} \varphi\right) v\right\}\right\}^{2} \\
\geq & -2\left(\square_{a} v\right) \tau \psi v-4\left(\sum_{j=1}^{n+1} A_{j} \partial_{j}^{2} v\right)\left(\sum_{j=1}^{n+1} A_{j} \tau \partial_{j} \varphi \partial_{j} v\right) \\
& +\left(2 \tau \psi \sum_{j=1}^{n+1} A_{j}\left(\tau^{2}\left(\partial_{j} \varphi\right)^{2}-\tau \partial_{j}^{2} \varphi\right) v^{2}-\tau^{2} \psi^{2} v^{2}\right)  \tag{4.18}\\
& -4\left(\sum_{j=1}^{n+1} A_{j} \tau \partial_{j} \varphi \partial_{j} v\right)\left(\sum_{k=1}^{n+1} A_{k}\left(\tau^{2}\left(\partial_{k} \varphi\right)^{2}-\tau \partial_{k}^{2} \varphi\right) v\right) \\
\equiv & S_{1}+S_{2}+S_{3}+S_{4},
\end{align*}
$$

where we have used an elementary inequality:

$$
\begin{aligned}
(B-C+D)^{2}= & (B+D-\tau \psi v)^{2}+C^{2}+2 B \tau \psi v \\
& -2 B C+2 \tau \psi v D-\tau^{2} \psi^{2} v^{2}-2 C D \\
\geq & 2 B \tau \psi v-2 B C+2 \tau \psi v D-\tau^{2} \psi^{2} v^{2}-2 C D
\end{aligned}
$$

with the obvious choice of $B, C, D$ and $S_{1}, S_{2}, S_{3}, S_{4}$.
Henceforth the sums are over $j, k=1, \ldots, n+1$, unless specified.
Multiplying the terms in $S_{2}$, we yield

$$
\begin{aligned}
S_{2}= & -4 \tau \sum A_{j} A_{k} \partial_{k} \varphi \partial_{j}^{2} v \partial_{k} v \\
= & -4 \tau \sum \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi \partial_{j} v \partial_{k} v\right)+2 \tau \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\left(\partial_{j} v\right)^{2}\right) \\
& -2 \tau \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} v\right)^{2}+4 \tau \sum \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} v \partial_{k} v,
\end{aligned}
$$

where we have used $2 \partial_{j} v\left(\partial_{k} \partial_{j} v\right)=\partial_{k}\left(\left(\partial_{j} v\right)^{2}\right)$.
Substituting $v=w u$, observing that $\partial_{j} v=w\left(\partial_{j} u+\tau\left(\partial_{j} \varphi\right) u\right)$ and generically denoting by $R_{0}$ the terms containing the factor $u$, and by $R, R_{1}, R_{2}, R_{3}$, $\ldots$. the terms bounded by $C w^{2}\left(|\nabla u|^{2}+\tau^{2} u^{2}\right.$ ) (where a constant $C>0$ is independent of $\tau$ ), we have

$$
\begin{aligned}
S_{2}= & 2 \tau \sum \partial_{j}\left(w^{2}\left(A_{j} A_{k} \partial_{j} \varphi\left(\partial_{k} u\right)^{2}-2 A_{j} A_{k} \partial_{k} \varphi \partial_{j} u \partial_{k} u+R_{0}\right)\right) \\
& -2 \tau w^{2} \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} u\right)^{2}-2 \tau^{2} w^{2} \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi \partial_{j}\left(u^{2}\right) \\
& -2 \tau^{3} w^{2} \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right)^{2} u^{2} \\
& +4 \tau w^{2} \sum \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} u \partial_{k} u \\
& +2 \tau^{2} w^{2} \sum \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi \partial_{k}\left(u^{2}\right)+\partial_{k} \varphi \partial_{j}\left(u^{2}\right)\right) \\
& +4 \tau^{3} w^{2} \sum \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi\left(\partial_{k} \varphi\right) u^{2} .
\end{aligned}
$$

We note

$$
\begin{aligned}
\tau^{2} w^{2} & \sum \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi \partial_{j}\left(u^{2}\right) \\
= & \sum \partial_{j}\left(\tau^{2} w^{2} \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right) u^{2}\right)-\tau^{2} \partial_{j}\left(w^{2}\right) \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right) u^{2} \\
& -\tau^{2} w^{2} \partial_{j} \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j}^{2} \varphi\right) u^{2} \\
= & -2 \tau^{3} w^{2}\left(\partial_{j} \varphi\right)^{2} \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right) u^{2}+\sum \partial_{k} R_{1}+R_{2} \\
= & \sum \partial_{k} R_{0}-2 \tau^{3} w^{2} \sum\left(\partial_{j} \varphi\right)^{2} \partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right) u^{2}+R .
\end{aligned}
$$

Moreover using similar relations for the terms with $\partial_{k}\left(u^{2}\right), \partial_{j}\left(u^{2}\right)$ and collecting
the terms with the factors $\tau w^{2}$ and $\tau^{3} w^{2}$, we have

$$
\begin{align*}
S_{2}= & 2 \tau \sum \partial_{j}\left(w^{2}\left(A_{j} A_{k} \partial_{j} \varphi\left(\partial_{k} u\right)^{2}-2 A_{j} A_{k} \partial_{k} \varphi \partial_{j} u \partial_{k} u\right)+R_{0}\right) \\
& +2 \tau w^{2} \sum\left\{2 \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} u \partial_{k} u-\partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} u\right)^{2}\right\}  \tag{4.19}\\
& +2 \tau^{3} w^{2} \sum\left\{\partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right)^{2}\right. \\
& \left.-2 \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi \partial_{k} \varphi\right\} u^{2}+R+\sum \partial_{j} R_{3} .
\end{align*}
$$

Similarly

$$
\begin{align*}
S_{3} & =2 \tau \psi w^{2} \sum A_{j}\left(\tau^{2}\left(\partial_{j} \varphi\right)^{2}-\tau \partial_{j}^{2} \varphi\right) u^{2}-\tau^{2} w^{2} \varphi^{2} u^{2} \\
& =2 \tau^{3} w^{2} \psi \sum A_{j}\left(\partial_{j} \varphi\right)^{2} u^{2}+R_{4} \tag{4.20}
\end{align*}
$$

and

$$
\begin{align*}
S_{4} & =-2 \tau \sum A_{j} A_{k} \partial_{j} \varphi\left(\tau^{2}\left(\partial_{k} \varphi\right)^{2}-\tau \partial_{k}^{2} \varphi\right) \partial_{j}\left(v^{2}\right) \\
& =2 \tau^{3} w^{2} \sum \partial_{j}\left(A_{j} A_{k} \partial_{j} \varphi\left(\partial_{k} \varphi\right)^{2}\right) u^{2}+\sum \partial_{j} R_{5}+R_{6} . \tag{4.21}
\end{align*}
$$

Noting $\partial_{j} v=\left(\partial_{j} u+\tau\left(\partial_{j} \varphi\right) u\right) w$, we see

$$
-2 \tau \sum A_{j} \psi\left(\partial_{j} v\right)^{2}=-2 \tau w^{2} \sum A_{j} \psi\left(\partial_{j} u+\tau\left(\partial_{j} \varphi\right) u\right)^{2}
$$

and

$$
\begin{aligned}
& -4 \tau^{2} w^{2} A_{j} \psi\left(\partial_{j} u\right)\left(\partial_{j} \varphi\right) u=-2 \tau^{2} w^{2}\left(A_{j} \psi\right)\left(\partial_{j}\left(u^{2}\right)\right)\left(\partial_{j} \varphi\right) \\
& \quad=-2 \tau^{2} \partial_{j}\left(w^{2}\left(A_{j} \psi \partial_{j} \varphi\right) u^{2}\right)+2 \tau^{2} \partial_{j}\left(w^{2}\right) A_{j} \psi\left(\partial_{j} \varphi\right) u^{2}+2 \tau^{2} w^{2} \partial_{j}\left(A_{j} \psi \partial_{j} \varphi\right) u^{2},
\end{aligned}
$$

we have

$$
\begin{align*}
S_{1} & =-2\left(\square_{a} v\right) \tau \psi v \\
& =2 \tau \sum \partial_{j}\left(A_{j} \psi\left(\partial_{j} v\right) v\right)-\tau \sum \partial_{j}\left(A_{j} \psi\right) \partial_{j}\left(v^{2}\right)-2 \tau \sum A_{j} \psi\left(\partial_{j} v\right)^{2}  \tag{4.22}\\
& =\sum \partial_{j} R_{0}+R-2 \tau w^{2} \sum A_{j} \psi\left(\partial_{j} u\right)^{2}+2 \tau^{3} w^{2} \sum A_{j} \psi\left(\partial_{j} \varphi\right)^{2} u^{2} .
\end{align*}
$$

By (4.19) through (4.22), we have

$$
\begin{align*}
w^{2}\left(\square_{a} u\right)^{2} \geq & 2 \tau \sum \partial_{j}\left(w^{2}\left(A_{j} A_{k}\left(\partial_{j} \varphi\left(\partial_{k} u\right)^{2}-2 \partial_{k} \varphi \partial_{j} u \partial_{k} u\right)\right)+R_{0}\right) \\
& +2 \tau w^{2} \sum\left\{2 \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} u \partial_{k} u\right. \\
& \left.-\partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} u\right)^{2}-A_{j} \psi\left(\partial_{j} u\right)^{2}\right\}  \tag{4.23}\\
& +2 \tau^{3} w^{2} \sum\left\{2 A_{j} \psi\left(\partial_{j} \varphi\right)^{2}+\partial_{k}\left(A_{j} A_{k} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right)^{2}\right. \\
& \left.-2 \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi \partial_{k} \varphi+\partial_{j}\left(A_{j} A_{k} \partial_{j} \varphi\left(\partial_{k} \varphi\right)^{2}\right)\right\} u^{2} \\
& +R+\sum \partial_{j} R_{0} .
\end{align*}
$$

Now we consider the integral of $w^{2}\left(\square_{a} u\right)^{2}$ over $Q(\delta)$. First we estimate the integral of the divergence terms in the right side of (4.23). Since $R_{0}$ contains $u$ as factor, the divergence theorem yields

$$
\int_{Q(\delta)} \partial_{j} R_{0} d x d t=0
$$

by $u_{\mid \partial Q(\delta)}=0$. Moreover, since $u=0$ on $\partial Q(\delta)$, as the surface integrals, we have only

$$
2 \tau \int_{\partial Q(\delta) \cap(\Gamma \times(0, T))} \sum w^{2} A_{j} A_{k}\left\{\partial_{j} \varphi\left(\partial_{k} u\right)^{2}-2 \partial_{k} \varphi \partial_{j} u \partial_{k} u\right\} \nu_{j} d \sigma .
$$

When $k=n+1$, the integral is zero because $\partial_{n+1} u=0$ on $\Gamma \times(0, T)$ and when $j=n+1$, it is zero because $\nu_{n+1}=0$ there. Due to the condition $u=0$ on $\Gamma \times(0, T)$, we have

$$
\partial_{j} u=\nu_{j} \partial_{\nu} u
$$

Hence the integral over $\partial Q(\delta) \cap(\Gamma \times(0, T))$ is

$$
\begin{aligned}
2 \tau \int_{\partial Q(\delta) \cap(\Gamma \times(0, T))} w^{2} & \sum\left(\left(\partial_{j} \varphi\right) \nu_{k}^{2}-2\left(\partial_{k} \varphi\right) \nu_{j} \nu_{k}\right)\left(\partial_{\nu} u\right)^{2} \nu_{j} d \sigma \\
& =-2 \tau \int_{\partial Q(\delta) \cap(\Gamma \times(0, T))} w^{2} \sum^{n}\left(\partial_{j} \varphi\right) \nu_{j}\left(\partial_{\nu} u\right)^{2} d \sigma
\end{aligned}
$$

because $\sum^{n} \nu_{j}^{2}=1$. According to our definition of $\varphi$, we have

$$
\sum^{n} \partial_{j} \varphi \cdot \nu_{j}=\sum^{n} 2\left(x-x_{0}\right)_{j} \cdot \nu_{j} \leq 0
$$

on $\Gamma_{-}$. Summing up, we obtain
[the surface integrals from the right side of (4.23)]

$$
\geq-2 \tau C \int_{\partial Q(\delta) \cap\left(\Gamma_{+} \times(0, T)\right)} w^{2}\left(\partial_{\nu} u\right)^{2} d \sigma
$$

with some constant $C>0$. Henceforth, generically by $C$, we denote a positive constant independent of $\tau$.

We break the factor of $2 \tau w^{2}$ in (4.23) into the sums over $j, k=1, \ldots, n$; $j=1, \ldots, n, k=n+1 ; k=1, \ldots, n, j=n+1 ; j=k=n+1$ to conclude that
this factor is

$$
\begin{aligned}
\sum^{n} & \left\{2 \partial_{j}^{2} \varphi\left(\partial_{j} u\right)^{2}-\partial_{k}^{2} \varphi\left(\partial_{j} u\right)^{2}-\psi\left(\partial_{j} u\right)^{2}\right\} \\
& +\sum^{n}\left\{2\left(\partial_{j}(-a) \partial_{n+1} \varphi\right) \partial_{j} u \partial_{n+1} u-\partial_{n+1}\left((-a) \partial_{n+1} \varphi\right)\left(\partial_{j} u\right)^{2}\right\} \\
& +\sum^{n}\left\{2 \partial_{n+1}\left(-a \partial_{k} \varphi\right) \partial_{n+1} u \partial_{k} u+\partial_{k}\left(a \partial_{k} \varphi\right)\left(\partial_{n+1} u\right)^{2}\right\}+a \psi\left(\partial_{n+1} u\right)^{2} \\
& +\partial_{n+1}\left(a^{2} \partial_{n+1} \varphi\right)\left(\partial_{n+1} u\right)^{2} \\
= & (4-2 n-\psi)\left|\nabla_{x} u\right|^{2}+\sum^{n}\left\{4 \beta\left(t-\frac{T}{2}\right) \partial_{j} a \partial_{j} u \partial_{t} u-2 \beta a\left(\partial_{j} u\right)^{2}\right\} \\
& +\sum 2\left(\partial_{j} a\left(x-x_{0}\right)_{j}+n a\right)\left(\partial_{t} u\right)^{2}+a \psi\left(\partial_{t} u\right)^{2}-2 \beta a^{2}\left(\partial_{t} u^{2}\right)
\end{aligned}
$$

where we have used definition (4.14) of $\varphi$, the equality $\partial_{n+1}=\partial_{t}$ and time independence of $a$. By using the Schwarz inequality for the scalar product $\nabla_{x} a \cdot \nabla_{x} u$ and then the inequality $2\left|\nabla_{x} u\right|\left|\partial_{t} u\right| \leq\left|\nabla_{x} u\right|^{2}+\left|\partial_{t} u\right|^{2}$, we conclude that the last sum is not less than

$$
\begin{aligned}
\{4 & \left.-2 n-\psi-2 \beta\left|\nabla_{x} a\right|\left|t-\frac{T}{2}\right|-2 \beta a\right\}\left|\nabla_{x} u\right|^{2} \\
& +\left\{-2 \beta\left|t-\frac{T}{2}\right|\left|\nabla_{x} a\right|+2\left(\nabla_{x} a \cdot\left(x-x_{0}\right)+n a\right)+a \psi-2 \beta a^{2}\right\}\left(\partial_{t} u\right)^{2} \\
\geq & c\left|\nabla_{x} u\right|^{2} \\
& +\left\{-2 \beta|\nabla a|\left|t-\frac{T}{2}\right|(1+a)+2 \nabla a \cdot\left(x-x_{0}\right)+(4-\varepsilon) a-4 \beta a^{2}\right\}\left(\partial_{t} u\right)^{2}
\end{aligned}
$$

due to choice (4.16) of $\psi$. Summing up, we can claim that for some small $\varepsilon$, in view of (4.1), the factor of $\tau w^{2}$ in (4.23) is not less than

$$
c\left|\nabla_{x} u\right|^{2}+c\left|\partial_{t} u\right|^{2}
$$

with some constant $c>0$, provided that

$$
\beta<\frac{2\left(\nabla a \cdot\left(x-x_{0}\right)\right)+4 a}{T|\nabla a|(1+a)+4 a^{2}} \quad \text { on } \quad \bar{Q} .
$$

Similarly the coefficient of $\tau^{3} w^{2}$ in (4.23) is
[the coefficient of $2 \tau^{3} w^{2} u^{2}$ ]

$$
\begin{aligned}
= & \sum\left\{2 A_{j} \psi\left(\partial_{j} \varphi\right)^{2}+\partial_{k}\left(A_{k} A_{j} \partial_{k} \varphi\right)\left(\partial_{j} \varphi\right)^{2}\right. \\
& \left.-2 \partial_{j}\left(A_{j} A_{k} \partial_{k} \varphi\right) \partial_{j} \varphi \partial_{k} \varphi+\partial_{j}\left(A_{j} A_{k} \partial_{j} \varphi\left(\partial_{k} \varphi\right)^{2}\right)\right\} \\
= & 8\left\{(\psi+2 n+2 a \beta)\left|x-x_{0}\right|^{2}-\left(a \psi \beta^{2}+2 n a \beta^{2}+2 a^{2} \beta^{3}\right)\left|t-\frac{T}{2}\right|^{2}\right\} \\
= & 8\left\{\left(4-\varepsilon-2 \beta|\nabla a|\left|t-\frac{T}{2}\right|\right)\left|x-x_{0}\right|^{2}\right. \\
& \left.-\left(\beta^{2}(4-\varepsilon) a-2 a \beta^{3}|\nabla a|\left|t-\frac{T}{2}\right|\right)\left|t-\frac{T}{2}\right|^{2}\right\} \\
> & 8\left\{2 \beta | t - \frac { T } { 2 } | ^ { 2 } \left(a \beta^{2}|\nabla a|\left|t-\frac{T}{2}\right|\right.\right. \\
& \left.-\left(|\nabla a|\left|t-\frac{T}{2}\right|+a\left(2-\frac{\varepsilon}{2}\right)\right) \beta+\left(2-\frac{\varepsilon}{2}\right)\right) \\
& \left.+\left(4-\varepsilon-2 \beta|\nabla a|\left|t-\frac{T}{2}\right|\right) \delta^{2}\right\} .
\end{aligned}
$$

Here we have used $\varphi(t, x)>\delta^{2}$ in $Q(\delta)$. The first and the third conditions in (4.2) imply that

$$
a \beta^{2}|\nabla a|\left|t-\frac{T}{2}\right|-\left(|\nabla a|\left|t-\frac{T}{2}\right|+2 a\right) \beta+2>0 \quad \text { on } \quad \bar{Q} .
$$

Then for sufficiently small $\varepsilon>0$, we obtain

$$
\begin{aligned}
a \beta^{2}|\nabla a|\left|t-\frac{T}{2}\right| & -\left(|\nabla a|\left|t-\frac{T}{2}\right|+a\left(2-\frac{\varepsilon}{2}\right)\right) \beta \\
& +\left(2-\frac{\varepsilon}{2}\right) \geq 0 \quad \text { on } \quad \bar{Q} .
\end{aligned}
$$

Therefore, under the first condition in (4.2),

$$
\begin{aligned}
{\left[\text { the coefficient of } 2 \tau^{3} w^{2} u^{2}\right] } & \geq 8\left(4-\varepsilon-2 \beta|\nabla a|\left|t-\frac{T}{2}\right|\right) \delta^{2} \\
& \geq 8(4-\varepsilon-T \beta|\nabla a|) \delta^{2} .
\end{aligned}
$$

Hence we choose a sufficiently small $\varepsilon>0$, so that we can complete the proof of Theorem 4.1.

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